

Bernoulli **21**(4), 2015, 2336–2350

DOI: [10.3150/14-BEJ646](https://doi.org/10.3150/14-BEJ646)

Density convergence in the Breuer–Major theorem for Gaussian stationary sequences

YAOZHONG HU^{1,*}, DAVID NUALART^{1,**}, SAMY TINDEL² and FANGJUN XU³

¹*Department of Mathematics, University of Kansas, Lawrence, KS 66045, USA.*

E-mail: ^{*}yhu@ku.edu; ^{**}nualart@ku.edu

²*Institut Élie Cartan, Université de Lorraine, 54506 Vandoeuvre-lès-Nancy, France.*

E-mail: samy.tindel@univ-lorraine.fr

³*School of Finance and Statistics, East China Normal University, Shanghai 200241, China.*

E-mail: fangjunxu@gmail.com

Consider a Gaussian stationary sequence with unit variance $X = \{X_k; k \in \mathbb{N} \cup \{0\}\}$. Assume that the central limit theorem holds for a weighted sum of the form $V_n = n^{-1/2} \sum_{k=0}^{n-1} f(X_k)$, where f designates a finite sum of Hermite polynomials. Then we prove that the uniform convergence of the density of V_n towards the standard Gaussian density also holds true, under a mild additional assumption involving the causal representation of X .

Keywords: Breuer–Major theorem; density convergence; Gaussian stationary sequences; Malliavin calculus; moving average representation

1. Introduction

Let $X = \{X_k; k \in \mathbb{N} \cup \{0\}\}$ be a centered Gaussian stationary sequence with unit variance. For all $v \in \mathbb{Z}$, we set $\rho(v) = \mathbf{E}[X_0 X_{|v|}]$. Therefore, $\rho(0) = 1$ and $|\rho(v)| \leq 1$ for all v . Let γ be the standard Gaussian probability measure and $f \in L^2(\gamma)$ be a fixed deterministic function such that $\mathbf{E}[f(X_1)] = 0$. We expand f in the orthonormal basis of Hermite polynomials $\{H_k; k \geq 0\}$, which are more specifically defined in Section 2.2. In particular, if f has Hermite rank $d \geq 1$, it admits the following Hermite expansion:

$$f(x) = \sum_{j=d}^{\infty} a_j H_j(x),$$

with $a_d \neq 0$. Define $V_n = \frac{1}{\sqrt{n}} \sum_{k=0}^{n-1} f(X_k)$. Then the celebrated Breuer–Major theorem (see [3] or Theorem 7.2.4 in [14]) can be written as follows:

This is an electronic reprint of the original article published by the ISI/BS in *Bernoulli*, 2015, Vol. 21, No. 4, 2336–2350. This reprint differs from the original in pagination and typographic detail.

Theorem 1.1. *Suppose that $\sum_{v \in \mathbb{Z}} |\rho(v)|^d < \infty$ and set $\sigma^2 = \sum_{j=d}^{\infty} j! a_j^2 \sum_{v \in \mathbb{Z}} \rho(v)^j$, which is assumed to be in $(0, \infty)$. Then the convergence*

$$V_n \xrightarrow{Law} \mathcal{N}(0, \sigma^2) \quad (1)$$

holds true as n tends to infinity.

We shall be in fact interested in a particular case of Theorem 1.1 for finite linear combinations of Hermite polynomials, which is stated here for convenience.

Corollary 1.2. *Consider $2 \leq d \leq q < \infty$ and a family of real numbers $\{a_j; j = d, \dots, q\}$. Let H_j be the j th order Hermite polynomial, and assume that $\sigma^2 \in (0, \infty)$, where $\sigma^2 \equiv \sum_{j=d}^q j! a_j^2 \sum_{v \in \mathbb{Z}} \rho(v)^j$. Set*

$$V_n^{d,q} = \frac{1}{\sqrt{n}} \sum_{k=0}^{n-1} \sum_{j=d}^q a_j H_j(X_k). \quad (2)$$

Then $V_n^{d,q} \xrightarrow{Law} \mathcal{N}(0, \sigma^2)$ as n tends to infinity. In particular, we have

$$\lim_{n \rightarrow \infty} \mathbf{E}[(V_n^{d,q})^4] = 3\sigma^4. \quad (3)$$

Remark 1.3. The relation between Gaussian convergence in law for sequences in a fixed Wiener chaos and behavior of the fourth moment has been extensively studied since the seminal paper [18]. We will need only a small part of the information available on the topic, such as relation (3).

Due to its importance, Breuer–Major theorem has been extended and refined in several directions. Important generalizations can be found in Arcones [1] (multidimensional case), Chambers and Slud [6] and Giraitis and Surgailis [7]. A proof of Theorem 1.1 using a combination of Stein’s method with Malliavin calculus was given by Nourdin, Peccati and Podolskij in [15], where one can find explicit bounds in the total variation and Wasserstein distances. We refer the reader to the monograph by Nourdin and Peccati [13] for a more detailed account on this topic.

We shall mainly be concerned here by convergences of densities, and here again the relationship between fourth moment behavior and various type of convergences of random variables in a fixed Wiener chaos have been thoroughly studied in the recent past. The interested reader is referred to [14] for further details, but we will use here the following recent criterion (see [12], Corollary 1.2, and [11], Corollary 4.6).

Theorem 1.4. *Let $\{F_n; n \in \mathbb{N}\}$ be a sequence of random variables belonging to a fixed chaos \mathcal{H}_q with $q \geq 2$. Suppose $\mathbf{E}[F_n^2] = 1$ and $\lim_{n \rightarrow \infty} \mathbf{E}[F_n^4] = 3$. Let p_{F_n} be the density of the random variable F_n and let $\phi(x) = (2\pi)^{-1/2} \exp(-|x|^2/2)$ be the density of the standard Gaussian distribution on \mathbb{R} .*

(i) Suppose that for some $\varepsilon > 0$,

$$\sup_n \mathbf{E}[\|DF_n\|^{-4-\varepsilon}] < \infty.$$

Then, there exists a constant c such that for all $n \geq 1$,

$$\sup_{x \in \mathbb{R}} |p_{F_n}(x) - \phi(x)| \leq c \sqrt{\mathbf{E}[F_n^4] - 3}.$$

(ii) Suppose that for all $p \geq 1$,

$$\sup_n \mathbf{E}[\|DF_n\|^{-p}] < \infty.$$

Then, for any $m \geq 0$, there exists a constant c_m such that for all $n \geq 1$,

$$\sup_{x \in \mathbb{R}} |p_{F_n}^{(m)}(x) - \phi^{(m)}(x)| \leq c_m \sqrt{\mathbf{E}[F_n^4] - 3}.$$

The goal of the current paper is to apply the criterion given by Theorem 1.4 in order to get convergence of density in the landmark of Breuer–Major theorem. In order to do this, we need a uniform estimate on the negative moments of the Malliavin derivative of the sequence, and this is the contents of our main result.

Theorem 1.5. *Let X be a Gaussian stationary sequence whose spectral density f_ρ satisfies $\log(f_\rho) \in L^1([-\pi, \pi])$ (see Hypothesis 2.1 and the examples in the next section). Let $V_n^{d,q}$ be the random variable defined by (2), and assume the hypothesis of Corollary 1.2 to be satisfied. Then for any $p \geq 1$, there exists n_0 such that*

$$\sup_{n \geq n_0} \mathbf{E}[\|DV_n^{d,q}\|^{-p}] < \infty. \quad (4)$$

In the case of a fixed Wiener chaos, we can obtain the following consequence.

Corollary 1.6. *Under the conditions of Theorem 1.5, if $q = d$, and we define $F_n = V_n^{d,d}/\sigma_n$, where $\sigma_n^2 = \mathbf{E}[(V_n^{d,d})^2]$, then, for all $m \geq 0$ there exists an n_0 (depending on m) such that*

$$\sup_{n \geq n_0} \sup_{x \in \mathbb{R}} |p_{F_n}^{(m)}(x) - \phi^{(m)}(x)| \leq c_m \sqrt{\mathbf{E}[F_n^4] - 3}.$$

In the case $q \neq d$, Theorem 1.4 cannot be applied. In the reference [11] one can find results on the uniform convergence of density for general random variables similar to those stated in Theorem 1.4, but they require a uniform lower bound for the negative moments of the random variable $|\langle DF_n, -DL^{-1}F_n \rangle_{\mathfrak{H}}|$, and our approach does not seem to work in this case because it is not clear how to express $\langle DF_n, -DL^{-1}F_n \rangle_{\mathfrak{H}}$ as a sum of squares. Nevertheless, condition (4) allows us to derive the uniform convergence of the

densities and their derivatives from a general result proved below (see Proposition 2.6) although in this case we have no information about the rate of convergence.

Corollary 1.7. *Under the conditions of Theorem 1.5, if we define $F_n = V_n^{d,q}/\sigma_n$, where $\sigma_n^2 = \mathbf{E}[(V_n^{d,q})^2]$, then, for all $m \geq 0$ we have*

$$\lim_{n \rightarrow \infty} \sup_{x \in \mathbb{R}} |p_{F_n}^{(m)}(x) - \phi^{(m)}(x)| = 0.$$

Notice that a particular case of Theorem 1.5 has been established in [12], for $q = 2$ and $X_k = B_{k+1} - B_k$ for a fractional Brownian motion B with Hurst parameter $H \in (0, 1)$. The proof of the existence of negative moments for $\|DF_n\|$ there is based on the Volterra representation of B , which leads to long computations. In comparison our current Theorem 1.5 is more general, since it is valid for a wide class of Gaussian stationary sequences. Its proof is also significantly simplified. These are achieved by the introduction of two new ingredients in the proof, namely:

- A general formula to compute conditional expectations for random variables of the form $H_q(X_k)$.
- Related to the previous item, we heavily resort to the causal representation of X_k , which is particularly convenient in order to compute conditional expectations.

Here is how our paper is structured: we give some preliminary results concerning Gaussian stationary sequences and related Malliavin calculus in Section 2. We then prove our main Theorem 1.5 in Section 3.

2. Preliminaries

This section is devoted to some preliminaries on causal or moving average representations for Gaussian stationary sequences, as well as Malliavin calculus tools which will be used in the sequel.

2.1. Moving average representation

The classical results on time series presented here are borrowed from [2, 4], to which we refer for further details. Start from our Gaussian stationary sequence $\{X_k; k \in \mathbb{N} \cup \{0\}\}$ with covariance function ρ . We will work under the following assumptions:

Hypothesis 2.1. *We suppose that ρ admits a spectral density f_ρ such that $\log(f_\rho) \in L^1([-\pi, \pi])$.*

Condition $\log(f_\rho) \in L^1([-\pi, \pi])$ is referred to as *purely nondeterministic* property in the literature. The interest of dealing with purely nondeterministic sequences is that they admit a so-called causal representation which is particularly convenient for conditional

expectation computations. Let us state a result in this direction, which is taken from [4], Chapter 5.

Proposition 2.2. *Let X be a Gaussian stationary sequence satisfying Hypothesis 2.1. Then for each $k \in \mathbb{N} \cup \{0\}$ the random variable X_k can be decomposed as*

$$X_k = \sum_{j \geq 0} \psi_j w_{k-j}, \quad (5)$$

where $(w_k)_{k \in \mathbb{Z}}$ is a discrete Gaussian white noise and the coefficients ψ_j are deterministic. We can always choose the white noise and the coefficients in such a way that $\psi_0 > 0$.

Proof. The existence of the causal moving average representation is a classical result that can be found, for instance, in [4], Theorem 7.5.2. \square

Notice that from the expansion (5) we easily obtain:

$$\rho(k_1 - k_2) = \sum_{l=-\infty}^{k_1 \wedge k_2} \psi_{k_1-l} \psi_{k_2-l}, \quad \text{for all } k_1, k_2 \in \mathbb{N}, \quad (6)$$

and this relation will be used in the proof of our main theorem.

Let us now turn to examples for which our standing assumptions of Hypothesis 2.1 are met. The following proposition provides two typical and classical cases for which a spectral density exists and satisfies some integrability properties.

Proposition 2.3. *Let ρ be the covariance function of X . We have the following statements.*

(i) *If $\rho \in \ell^1$, then the spectral density f_ρ exists and is a nonnegative bounded function defined on $[-\pi, \pi]$.*

(ii) *Suppose that $\{\rho(k)k^\alpha, k > 0\}$ is positive and it is normalized slowly varying at infinity for some $\alpha \in (0, 1)$. That is, for every $\delta > 0$, for sufficiently large k , $\rho(k)k^{\alpha+\delta}$ is increasing and $\rho(k)k^{\alpha-\delta}$ is decreasing. Then the spectral density exists and satisfies $\lim_{\lambda \rightarrow 0} |\lambda|^{1-\alpha} f_\rho(\lambda) = c_f$ for some constant $c_f > 0$ (see [9]).*

We now give two specific and important examples which satisfy Hypothesis 2.1.

Example 2.4. The so-called autoregressive fractionally integrated moving-average (ARFIMA) processes are introduced in [8] and [10]. Denote by B the one lag backward operator ($BX_k = X_{k-1}$). Let $\phi(z)$ and $\theta(z)$ be two polynomials which have no common zeros and such that the zeros of ϕ lie outside the closed unit disk $\{z, |z| \leq 1\}$. Suppose that X_k is given by

$$\phi(B)X_k = (\text{Id} - B)^{-d} \theta(B)w_k, \quad (7)$$

where $-1 < d < 1/2$, and where the operator $(\text{Id} - B)^{-d}$ is defined by:

$$(\text{Id} - B)^{-d} = \sum_{j=1}^{\infty} \eta_j B^j \quad \text{with } \eta_j = \frac{\Gamma(d+j)}{\Gamma(j+1)\Gamma(d)}.$$

Also notice that in (7) the sequence $(w_k)_{k \in \mathbb{Z}}$ is a discrete Gaussian white noise. It is well-known (see [19], Theorem 3.4 and equation (3.19)) that under the above conditions, $\{X_k, k \in \mathbb{N}\}$ admits a spectral density whose exact expression is:

$$f(\lambda) = \frac{1}{2\pi} \left[2 \sin \frac{\lambda}{2} \right]^{-2d} \frac{|\theta(e^{-i\lambda})|^2}{|\phi(e^{-i\lambda})|^2}.$$

It is thus readily checked that Hypothesis 2.1 is satisfied, and hence X_k has a causal representation.

Example 2.5. Our second example is the fractional Gaussian noise. Let $\{B_t, t \geq 0\}$ be a fractional Brownian motion of Hurst parameter $H \in (0, 1)$. Then $\{X_k = B_{k+1} - B_k, k \in \mathbb{N} \cup \{0\}\}$ is a stationary Gaussian process with correlation

$$\rho(k) = \frac{1}{2} [(k+1)^{2H} - 2k^{2H} + (k-1)^{2H}].$$

Its spectral density (see, e.g., [2], equation (2.17)) is:

$$f(\lambda) = \frac{1}{2\pi} \sum_{k=-\infty}^{\infty} \rho(|k|) e^{ik\lambda} = 2c_f (1 - \cos(\lambda)) \sum_{j=-\infty}^{\infty} |2\pi j + \lambda|^{-2H-1}, \quad \lambda \in [-\pi, \pi],$$

where $c_f = (2\pi)^{-1} \sin(\pi H) \Gamma(2H+1)$. If $H \leq 1/2$, it is clear that $\sum_{k=-\infty}^{\infty} |\rho(|k|)| < \infty$. This implies

$$\sup_{\lambda \in [-\pi, \pi]} |f(\lambda)| < \infty.$$

If $1/2 < H < 1$, then

$$0 \leq f(\lambda) \leq 2c_f (1 - \cos(\lambda)) |\lambda|^{-2H-1} + 2c_f \sum_{j \neq 0} |2\pi j + \lambda|^{-2H-1}, \quad \lambda \in [-\pi, \pi].$$

The first term is in L^1 since $H < 1$. When $j \neq 0$, $\int_{-\pi}^{\pi} |2\pi j + \lambda|^{-2H-1} d\lambda \leq C j^{-2H}$ for some positive constant C . Thus $\int_{-\pi}^{\pi} \sum_{j \neq 0} |2\pi j + \lambda|^{-2H-1} d\lambda < \infty$, owing to the fact that $H > 1/2$. Therefore, we have $f \in L^1$. Summarizing we have $f \in L^1$ for all $H \in (0, 1)$. This also implies $\log^+ f(\lambda) \in L^1$. To see $\log^- f(\lambda) \in L^1$, we notice that

$$f(\lambda) \geq 2c_f (1 - \cos(\lambda)) |\lambda|^{-2H-1}.$$

So $\log^- f(\lambda) \leq C + |\log[(1 - \cos(\lambda)) |\lambda|^{-2H-1}]|$ which is in L^1 . In conclusion, the sequence X satisfies Hypothesis 2.1.

2.2. Malliavin calculus

We start by briefly recalling some basic notation and results connected to Gaussian analysis and Malliavin calculus. The reader is referred to [14, 17] for details or missing proofs.

2.2.1. Wiener space and generalities

Let \mathfrak{H} be a real separable Hilbert space with inner product $\langle \cdot, \cdot \rangle_{\mathfrak{H}}$. The norm of \mathfrak{H} will be denoted by $\| \cdot \| = \| \cdot \|_{\mathfrak{H}}$. Recall that we call *isonormal Gaussian process* over \mathfrak{H} any centered Gaussian family $W = \{W(h) : h \in \mathfrak{H}\}$, defined on a probability space $(\Omega, \mathcal{F}, \mathbf{P})$ and such that $\mathbf{E}[W(h)W(g)] = \langle h, g \rangle_{\mathfrak{H}}$ for every $h, g \in \mathfrak{H}$.

In our application the underlying Gaussian family will be a discrete Gaussian white noise $(w_k)_{k \in \mathbb{Z}}$. The space \mathfrak{H} is given here by $\mathfrak{H} = \ell^2(\mathbb{Z})$ (the space of square integrable sequences indexed by \mathbb{Z}) equipped with its natural inner product. Set $\{\varepsilon^j; j \in \mathbb{Z}\}$ for the canonical basis of $\ell^2(\mathbb{Z})$, that is $\varepsilon_k^j = \delta_j(k)$. We thus identify w_j with $W(\varepsilon^j)$. Assume from now on that our underlying σ -algebra \mathcal{F} is generated by W .

For any integer $q \in \mathbb{N} \cup \{0\}$, we denote by \mathcal{H}_q the q th *Wiener chaos* of W . We recall that \mathcal{H}_0 is simply \mathbb{R} whereas, for any $q \geq 1$, \mathcal{H}_q is the closed linear subspace of $L^2(\Omega)$ generated by the family of random variables $\{H_q(W(h)), h \in \mathfrak{H}, \|h\|_{\mathfrak{H}} = 1\}$, with H_q the q th Hermite polynomial given by

$$H_q(x) = (-1)^q e^{x^2/2} \frac{d^q}{dx^q} (e^{-x^2/2}). \quad (8)$$

Let \mathcal{S} be the set of all cylindrical random variables of the form

$$F = g(W(h_1), \dots, W(h_n)),$$

where $n \geq 1$, $h_i \in \mathfrak{H}$, and g is infinitely differentiable such that all its partial derivatives have polynomial growth. The Malliavin derivative of F is the element of $L^2(\Omega; \mathfrak{H})$ defined by

$$DF = \sum_{i=1}^n \frac{\partial g}{\partial x_i} (W(h_1), \dots, W(h_n)) h_i.$$

By iteration, for every $m \geq 2$, we define the m th derivative $D^m F$. This is an element of $L^2(\Omega; \mathfrak{H}^{\odot m})$, where $\mathfrak{H}^{\odot m}$ designates the symmetric m th tensor product of \mathfrak{H} . For $m \geq 1$ and $p \geq 1$, $\mathbb{D}^{m,p}$ denote the closure of \mathcal{S} with respect to the norm $\| \cdot \|_{m,p}$ defined by

$$\|F\|_{m,p}^p = \mathbf{E}[|F|^p] + \sum_{j=1}^m \mathbf{E}[\|D^j F\|_{\mathfrak{H}^{\otimes j}}^p].$$

Set $\mathbb{D}^\infty = \bigcap_{m,p} \mathbb{D}^{m,p}$. One can then extend the definition of D^m to $\mathbb{D}^{m,p}$. When $m = 1$, one simply write D instead of D^1 . As a consequence of the hypercontractivity property of

the Ornstein–Uhlenbeck semigroup (see, e.g., [14], Theorem 2.7.2), all the $\|\cdot\|_{m,p}$ -norms are equivalent in any *finite* sum of Wiener chaoses.

Finally, let us recall that the Malliavin derivative D satisfies the following *chain rule*: if $\varphi: \mathbb{R}^n \rightarrow \mathbb{R}$ is in \mathcal{C}_b^1 (i.e., belongs to the set of continuously differentiable functions with a bounded derivative) and if $\{F_i\}_{i=1,\dots,n}$ is a vector of elements of $\mathbb{D}^{1,2}$, then $\varphi(F_1, \dots, F_n) \in \mathbb{D}^{1,2}$ and

$$D\varphi(F_1, \dots, F_n) = \sum_{i=1}^n \frac{\partial \varphi}{\partial x_i}(F_1, \dots, F_n) DF_i. \quad (9)$$

2.3. Convergence of densities

Suppose that F is a random variable in \mathbb{D}^∞ such that $\mathbf{E}[\|DF\|^{-p}] < \infty$ for all $p \geq 1$. Then, we know that F has an infinitely differentiable density and there are explicit formulas for the density and its derivatives (see [17], Proposition 2.1.5). Using this result, we can establish the following criteria for convergence of densities for random variables in a finite sum of Wiener chaoses.

Proposition 2.6. *Let $\{F_n; n \in \mathbb{N}\}$ be a sequence of random variables belonging to a finite sum of Wiener chaoses $\bigoplus_{k=1}^M \mathcal{H}_k$, which converges in law to a nonzero random variable F_∞ . Suppose that, for all $p \geq 1$*

$$\sup_n \mathbf{E}[\|DF_n\|^{-p}] < \infty. \quad (10)$$

Then, for all $m \geq 0$ the derivative $p_n^{(m)}$ of the density of F_n , converges uniformly and in $L^p(\mathbb{R})$ for all $p \geq 1$ to the corresponding derivative of the density of F_∞ .

Proof. First, notice that by condition (10), the random variable F_n has an infinitely differentiable density p_n , whose derivatives can be expressed by

$$p_n^{(m)}(x) = \mathbf{E}[\mathbf{1}_{\{F_n > x\}} G_n^{(m)}], \quad (11)$$

where the random variables $G_n^{(m)}$ are defined recursively by $G_n^{(0)} = \delta(\frac{DF_n}{\|DF_n\|_H^2})$ and

$$G_n^{(m)} = -\delta\left(\frac{G_n^{(m-1)} DF_n}{\|DF_n\|_H^2}\right),$$

for any $m \geq 1$. From this formula, it follows that the derivatives $p_n^{(m)}$ are uniformly bounded and also uniformly bounded in $L^p(\mathbb{R})$ for all $p \geq 1$. Indeed, by [16], Lemma 2.4, we have $\sup_n \mathbf{E}[|F_n|^r] < \infty$ for all $r \geq 1$. This uniform bound on the moments, together with the equivalence of the $\|\cdot\|_{m,p}$ norms in any finite sum of Wiener chaoses and

condition (10) imply that $\sup_n \|G_n^{(m)}\|_{L^p(\Omega)} = c_{m,p} < \infty$ for all $m \geq 0$. Then, we can write from (11)

$$\sup_n \sup_x |p_n^{(m)}(x)| \leq \sup_n \mathbf{E}[|G_n^{(m)}|] = c_{m,1} < \infty, \quad (12)$$

and, using the fact that $\mathbf{E}[G_n^{(m)}] = 0$, we get:

$$\begin{aligned} \sup_n |p_n^{(m)}(x)| &\leq \sup_n (\mathbf{P}(|F_n| > |x|) \mathbf{E}[|G_n^{(m)}|^2])^{1/2} \\ &\leq c_{m,2} \sup_n \sqrt{\mathbf{P}(|F_n| > |x|)} \\ &\leq c_{m,2} \sup_n \mathbf{E}[|F_n|^q]^{1/2} |x|^{-q/2} \leq c |x|^{-q/2}, \end{aligned} \quad (13)$$

for all $q \geq 1$ and for some constant c depending on q and m .

By [16], Theorem 3.1, the sequence F_n converges in total variation to F_∞ , that is, the densities p_n converge in $L^1(\mathbb{R})$ to the density p_∞ of F_∞ . The boundedness in $L^p(\mathbb{R})$ and the uniform bound of p_n imply that this convergence holds in $L^p(\mathbb{R})$ for any $p \geq 1$.

On the other hand, the estimates (12) and (13) imply that for any $m \geq 1$ and any $p \geq 1$, the sequence $p_n^{(m)}$ is uniformly bounded in $L^p(\mathbb{R})$. Therefore, for any $m \geq 1$ and any $p \geq 1$ the sequence $p_n^{(m)}$ is relatively compact in $L^p(\mathbb{R})$. Suppose that a subsequence $\{p_{n_k}^{(m)}, k \geq 1\}$ converges in $L^p(\mathbb{R})$ to some limit $\tilde{p}_\infty^{(m)}$. This limit must coincide with the m th derivative (in the distribution sense) of p_∞ , and, therefore, it is unique. This implies that for any $m \geq 1$ and any $p \geq 1$ the sequence $p_n^{(m)}$ converges in $L^p(\mathbb{R})$ to the m th derivative of p_∞ .

Finally, the uniform convergence is also easy to establish from the convergence of the densities in $L^p(\mathbb{R})$ for all $p \geq 1$. \square

2.4. A key lemma

Our future computations will heavily rely on an efficient way to compute conditional expectations. Towards this aim, we state here some general results. Let us start with a decomposition for Hermite polynomials.

Lemma 2.7. *For any $q \geq 1$, let H_q be the polynomial defined by relation (8). Consider $y, z \in \mathbb{R}$ and two real parameters a, b with $a^2 + b^2 = 1$. Then the following relation holds true:*

$$H_q(ay + bz) = \sum_{\ell=0}^q \binom{q}{\ell} a^{q-\ell} b^\ell H_{q-\ell}(y) H_\ell(z). \quad (14)$$

Proof. By the definition of the Hermite polynomials, we have

$$e^{aty-(at)^2/2} = \sum_{i=0}^{\infty} \frac{(at)^i}{i!} H_i(y) \quad \text{and} \quad e^{tbz-(bt)^2/2} = \sum_{j=0}^{\infty} \frac{(bt)^j}{j!} H_j(z). \quad (15)$$

In the same way, we also obtain

$$e^{t(ay+bz)-t^2/2} = \sum_{q=0}^{\infty} \frac{t^q}{q!} H_q(ay+bz). \quad (16)$$

Since $a^2 + b^2 = 1$, we obviously have $e^{aty-(at)^2/2} e^{tbz-(bt)^2/2} = e^{t(ay+bz)-t^2/2}$. Thus, multiplying the right-hand sides of both identities in (15) we recover the right-hand side of (16), namely:

$$\sum_{q=0}^{\infty} \frac{t^q}{q!} H_q(ay+bz) = \sum_{i=0}^{\infty} \frac{(at)^i}{i!} H_i(y) \sum_{j=0}^{\infty} \frac{(bt)^j}{j!} H_j(z),$$

which easily yields the desired identity (14). \square

With this preliminary result in hand, we are ready to state our result on conditional expectations.

Proposition 2.8. *Let Y and Z be two centered Gaussian random variables such that Y is measurable with respect to a σ -algebra $\mathcal{G} \subset \mathcal{F}$ and Z is independent of \mathcal{G} . Assume that $\mathbf{E}[Y^2] = \mathbf{E}[Z^2] = 1$. Then for any $q \geq 1$, and real parameters a, b such that $a^2 + b^2 = 1$, we have:*

$$\mathbf{E}[H_q(aY + bZ) | \mathcal{G}] = a^q H_q(Y). \quad (17)$$

Proof. Apply identity (14) in order to decompose $H_q(aY + bZ)$. Then identity (17) follows easily from the fact that Y is \mathcal{G} -measurable, Z is independent from \mathcal{G} and Hermite polynomials have 0 mean under a centered Gaussian measure except for $H_0 \equiv 1$. \square

2.5. Carbery–Wright inequality

In the proof of Theorem 1.5, we will make use of the following inequality due to Carbery and Wright [5], Theorem 8, which is recalled here for convenience.

Proposition 2.9. *Let $X = (X_1, \dots, X_n)$ be a Gaussian random vector in \mathbb{R}^n and $Q: \mathbb{R}^n \rightarrow \mathbb{R}$ a polynomial of degree at most m . Then there is a universal constant $c > 0$ such that:*

$$\mathbf{E}[|Q(X_1, \dots, X_n)|]^{1/m} \mathbf{P}(|Q(X_1, \dots, X_n)| \leq x) \leq cmx^{1/m}, \quad \text{for all } x > 0. \quad (18)$$

3. Proof of Theorem 1.5

In this section, we will prove our main result, which amounts to show the inequality (4). This will be done into several steps.

Step 1: Computation of the Malliavin norm. Recall that $V_n^{d,q}$ is defined by relation (2), and set for the moment $f = \sum_{j=d}^q a_j H_j$. Invoking relation (9), plus the fact that $Dw_j = \varepsilon^j$ with the notation of Section 2.2.1, we get:

$$DV_n^{d,q} = \frac{1}{\sqrt{n}} \sum_{k=0}^{n-1} f'(X_k) \left(\sum_{j \geq 0} \psi_j \varepsilon^{k-j} \right) = \frac{1}{\sqrt{n}} \sum_{l=-\infty}^{n-1} \left(\sum_{k=l+}^{n-1} \psi_{k-l} f'(X_k) \right) \varepsilon^l, \quad (19)$$

where $l^+ = \max\{l, 0\}$. Invoking relation (6), it is thus readily checked that:

$$\langle DV_n^{d,q}, DV_n^{d,q} \rangle_{\mathfrak{H}} = \frac{1}{n} \sum_{k_1, k_2=0}^{n-1} f'(X_{k_1}) \rho(k_1 - k_2) f'(X_{k_2}),$$

where we recall that ρ is the covariance function of the Gaussian stationary sequence $\{X_k; k \geq 0\}$. This is consistent with the expression found in [14], Chapter 5. However, in order to write the above expression as sum of some squares, we will start directly from expression (19). Since $\{\varepsilon^l; l \in \mathbb{Z}\}$ is an orthonormal basis of $\ell^2(\mathbb{Z})$ we obtain:

$$\langle DV_n^{d,q}, DV_n^{d,q} \rangle_{\mathfrak{H}} = \frac{1}{n} \sum_{\ell=-\infty}^{n-1} \left(\sum_{k=\ell+}^{n-1} \psi_{k-\ell} f'(X_k) \right)^2.$$

Rearranging terms (namely, change $k - \ell$ to k and then $n - \ell - 1$ to m), we end up with:

$$\begin{aligned} \langle DV_n^{d,q}, DV_n^{d,q} \rangle_{\mathfrak{H}} &\geq \frac{1}{n} \sum_{\ell=0}^{n-1} \left(\sum_{k=0}^{n-\ell-1} f'(X_{\ell+k}) \psi_k \right)^2 \\ &= \frac{1}{n} \sum_{m=0}^{n-1} \left(\sum_{k=0}^m f'(X_{n-1-(m-k)}) \psi_k \right)^2 \equiv A_n. \end{aligned}$$

As a last preliminary step we resort to the fact that $X = \{X_k; k \in \mathbb{N} \cup \{0\}\}$ is a Gaussian stationary sequence, which allows to assert that A_n is identical in law to B_n with

$$B_n := \frac{1}{n} \sum_{m=0}^{n-1} \left(\sum_{k=0}^m f'(X_{m-k}) \psi_k \right)^2 = \frac{1}{n} \sum_{m=0}^{n-1} \left(\sum_{k=0}^m f'(X_k) \psi_{m-k} \right)^2.$$

We will now bound the negative moments of B_n .

Step 2: Block decomposition. We now wish to apply the Carbery–Wright inequality (18) in order to get bounds for negative moments of B_n . However, relation (18) only

applies to moments of small order, and this is why we proceed to a decomposition of B_n into smaller blocks.

Fix thus an integer $N \geq 1$ and let $M = [n/N]$ be the integer part of n/N . Then $n \geq NM$ and as a consequence,

$$B_n = \frac{1}{n} \sum_{m=0}^{n-1} \left(\sum_{k=0}^m f'(X_k) \psi_{m-k} \right)^2 \geq \frac{1}{n} \sum_{i=0}^{N-1} \sum_{m=iM}^{(i+1)M-1} \left(\sum_{k=0}^m f'(X_k) \psi_{m-k} \right)^2.$$

For $i = 0, \dots, N-1$, define

$$B_n^i = \frac{1}{n} \sum_{m=iM}^{(i+1)M-1} \left(\sum_{k=0}^m f'(X_k) \psi_{m-k} \right)^2$$

so that $B_n \geq \sum_{i=0}^{N-1} B_n^i$. Then it is readily checked that:

$$(B_n)^{-p/2} \leq \prod_{i=0}^{N-1} (B_n^i)^{-p/(2N)}. \quad (20)$$

Recall once again the representation of the sequence X in (5), and denote by \mathcal{F}_k the filtration generated by $\{w_\ell : \ell < k\}$. Then starting from (20) we obtain:

$$\begin{aligned} \mathbf{E}[(B_n)^{-p/2}] &\leq \mathbf{E} \left[\prod_{i=0}^{N-1} (B_n^i)^{-p/(2N)} \right] \\ &= \mathbf{E} \left[\mathbf{E}[(B_n^{N-1})^{-p/(2N)} | \mathcal{F}_{(N-1)M}] \prod_{i=0}^{N-2} (B_n^i)^{-p/(2N)} \right]. \end{aligned} \quad (21)$$

Step 3: Application of Carbery–Wright. Let us go back to the particular situation of $f = \sum_{j=d}^q a_j H_j$, which means in particular that $f' = \sum_{j=d}^q j a_j H_{j-1}$. We are now in a position to apply a conditional version of inequality (18) to the block $(B_n^{N-1})^{-p/(2N)}$ in (21). First, we notice

$$\mathbf{E}[(B_n^{N-1})^{-p/(2N)} | \mathcal{F}_{(N-1)M}] \leq 1 + \frac{p}{2N} \int_0^1 \mathbf{P}(B_n^{N-1} \leq x | \mathcal{F}_{(N-1)M}) x^{-p/(2N)-1} dx. \quad (22)$$

Since B_n^{N-1} is a polynomial of order $m = 2(q-1)$, Carbery–Wright’s inequality (18) yields:

$$\mathbf{P}(B_n^{N-1} \leq x | \mathcal{F}_{(N-1)M}) \leq \frac{c x^{1/(2(q-1))}}{[\mathbf{E}(B_n^{N-1} | \mathcal{F}_{(N-1)M})]^{1/(2(q-1))}}. \quad (23)$$

Step 4: Estimates for the conditional expectation. We now estimate the conditional expectation $\mathbf{E}[B_n^{N-1}|\mathcal{F}_{(N-1)M}]$. We have:

$$\begin{aligned}\mathbf{E}[B_n^{N-1}|\mathcal{F}_{(N-1)M}] &= \frac{1}{n} \sum_{m=(N-1)M}^{NM-1} \mathbf{E} \left[\left(\sum_{k=0}^m f'(X_k) \psi_{m-k} \right)^2 \middle| \mathcal{F}_{(N-1)M} \right] \\ &\geq \frac{1}{n} \sum_{m=(N-1)M}^{NM-1} A_m,\end{aligned}\tag{24}$$

where we have set

$$A_m = \mathbf{Var} \left(\sum_{k=(N-1)M}^m f'(X_k) \psi_{m-k} \middle| \mathcal{F}_{(N-1)M} \right).$$

Furthermore, notice that

$$f'(X_k) = f' \left(\sum_{\ell=-\infty}^k \psi_{k-\ell} w_\ell \right) = f'(Y_k + Z_k),$$

where $Y_k = \sum_{i=-\infty}^{(N-1)M-1} \psi_{k-i} w_i$ is $\mathcal{F}_{(N-1)M}$ -measurable and $Z_k = \sum_{i=(N-1)M}^k \psi_{k-i} w_i$ is independent of $\mathcal{F}_{(N-1)M}$. Recalling that $f' = \sum_{j=d}^q j a_j H_{j-1}$, we can thus resort to Lemmas 2.7 and 2.8. This gives:

$$f'(X_k) - \mathbf{E}[f'(X_k)|\mathcal{F}_{(N-1)M}] = \sum_{j=d}^q \sum_{\ell=1}^{j-1} j a_j \binom{j-1}{\ell} \sigma_{Y_k}^{j-1-\ell} H_{j-1-\ell}(\tilde{Y}_k) \sigma_{Z_k}^\ell H_\ell(\tilde{Z}_k),$$

where $\sigma_{Y_k} = [\mathbf{Var}(Y_k)]^{1/2}$, $\sigma_{Z_k} = [\mathbf{Var}(Z_k)]^{1/2}$, $\tilde{Y}_k = Y_k/\sigma_{Y_k}$ and $\tilde{Z}_k = Z_k/\sigma_{Z_k}$. Therefore,

$$\begin{aligned}A_m &= \mathbf{E} \left[\left(\sum_{k=(N-1)M}^m \sum_{j=d}^q \sum_{\ell=1}^{j-1} a_{j,\ell,k} H_{j-1-\ell}(\tilde{Y}_k) H_\ell(\tilde{Z}_k) \psi_{m-k} \right)^2 \middle| \mathcal{F}_{(N-1)M} \right] \\ &= \mathbf{E} \left[\left(\sum_{\ell=1}^{q-1} \sum_{k=(N-1)M}^m \sum_{j=(\ell+1) \vee d}^q a_{j,\ell,k} H_{j-1-\ell}(\tilde{Y}_k) H_\ell(\tilde{Z}_k) \psi_{m-k} \right)^2 \middle| \mathcal{F}_{(N-1)M} \right],\end{aligned}$$

where we have set $a_{j,\ell,k} = j a_j \binom{j-1}{\ell} \sigma_{Y_k}^{j-1-\ell} \sigma_{Z_k}^\ell$.

Recall that the random variables \tilde{Y}_k are $\mathcal{F}_{(N-1)M}$ -measurable while the random variables \tilde{Z}_k are independent of $\mathcal{F}_{(N-1)M}$. By decorrelation properties of Hermite polynomials

we thus get:

$$A_m = \sum_{\ell=1}^{q-1} \mathbf{E} \left[\left(\sum_{k=(N-1)M}^m \sum_{j=(\ell+1) \vee d}^q a_{j,\ell,k} H_{j-1-\ell}(\tilde{Y}_k) H_\ell(\tilde{Z}_k) \psi_{m-k} \right)^2 \middle| \mathcal{F}_{(N-1)M} \right]$$

and we trivially lower bound this quantity by taking the term corresponding to $\ell = q-1$. In this situation, the sum $\sum_{j=(\ell+1) \vee d}^q$ is reduced to the term corresponding to $j = q$, and since $a_{q,q-1,k} = qa_q \sigma_{\tilde{Z}_k}^{q-1}$ we obtain:

$$\begin{aligned} A_m &\geq \mathbf{E} \left[\left(\sum_{k=(N-1)M}^m qa_q \sigma_{\tilde{Z}_k}^{q-1} H_{q-1}(\tilde{Z}_k) \psi_{m-k} \right)^2 \middle| \mathcal{F}_{(N-1)M} \right] \\ &= q^2 a_q^2 \mathbf{E} \left[\left(\sum_{k=(N-1)M}^m \sigma_{\tilde{Z}_k}^{q-1} H_{q-1}(\tilde{Z}_k) \psi_{m-k} \right)^2 \right]. \end{aligned}$$

We now invoke the identity $\mathbf{E}[H_p(\tilde{Z}_{k_1}) H_p(\tilde{Z}_{k_2})] = p! (\mathbf{E}[\tilde{Z}_{k_1} \tilde{Z}_{k_2}])^p$ in order to obtain

$$A_m \geq qq! a_q^2 \sum_{k_1, k_2=(N-1)M}^m \sigma_{\tilde{Z}_{k_1}}^{q-1} \sigma_{\tilde{Z}_{k_2}}^{q-1} \mathbf{E}[\tilde{Z}_{k_1} \tilde{Z}_{k_2}]^{q-1} \psi_{m-k_1} \psi_{m-k_2}.$$

Furthermore, similarly to (6), it is readily checked that:

$$\mathbf{E}[\tilde{Z}_{k_1} \tilde{Z}_{k_2}] = \frac{1}{\sigma_{\tilde{Z}_{k_1}} \sigma_{\tilde{Z}_{k_2}}} \sum_{i=(N-1)M}^{k_1 \wedge k_2} \psi_{k_1-i} \psi_{k_2-i},$$

and thus

$$\begin{aligned} A_m &\geq qq! a_q^2 \sum_{k_1, k_2=(N-1)M}^m \left(\sum_{i=(N-1)M}^{k_1 \wedge k_2} \psi_{k_1-i} \psi_{k_2-i} \right)^{q-1} \psi_{m-k_1} \psi_{m-k_2} \\ &= qq! a_q^2 \sum_{i_1, \dots, i_{q-1}=(N-1)M}^m \sum_{k_1, k_2=\max(i_1, \dots, i_{q-1})}^m \psi_{m-k_1} \psi_{m-k_2} \prod_{j=1}^{q-1} \psi_{k_1-i_j} \psi_{k_2-i_j} \\ &= qq! a_q^2 \sum_{i_1, \dots, i_{q-1}=(N-1)M}^m \left(\sum_{k=\max(i_1, \dots, i_{q-1})}^m \psi_{m-k} \prod_{j=1}^{q-1} \psi_{k-i_j} \right)^2. \end{aligned}$$

Here again, this sum of squares is trivially lower bounded by taking the term corresponding to $i_1 = \dots = i_{q-1} = m$, which yields:

$$A_m \geq c_{a,q,\psi} \quad \text{with } c_{a,q,\psi} \equiv qq! a_q^2 \psi_0^{2q} > 0. \quad (25)$$

Step 5: Conclusion. In the remainder of the proof the constants $c_{a,q,\psi,N}$ and so can change from line to line without further mention. Plugging relation (25) into (25) and recalling that N is a given integer whose exact value will be fixed below, we get:

$$\mathbf{E}[B_n^{N-1} | \mathcal{F}_{(N-1)M}] \geq \frac{M c_{a,q,\psi}}{n} \geq c_{a,q,\psi,N} > 0,$$

as long as N stays bounded. We then insert back this inequality into (22) and (23) in order to get:

$$\mathbf{P}(B_n^{N-1} \leq x | \mathcal{F}_{(N-1)M}) \leq 1 + \frac{p c_{a,q,\psi,N}}{2N} \int_0^1 x^{1/(2(q-1)) - p/(2N) - 1} dx = c_{a,q,\psi,N,p} < \infty,$$

where we have chosen N such that $\frac{p}{2N} < \frac{1}{2(q-1)}$. Iterating this bound into (21), we have thus obtained:

$$\mathbf{E}[(B_n)^{-p/2}] \leq c_{a,q,\psi,N,p}^N,$$

which is a finite quantity. Finally, recall from Step 1 that $\mathbf{E}[(B_n)^{-p/2}] \geq \mathbf{E}[\|DV_n^{d,q}\|_{\mathfrak{H}}^{-p}]$, which finishes the proof.

Acknowledgements

We would like to thank an anonymous referee for carefully reading this manuscript and making helpful remarks. The research of Y. Hu is partially supported by a grant from the Simons Foundation #209206 and by a General Research Fund of University of Kansas. The research of D. Nualart is supported by the NSF grant DMS1208625. The research of F. Xu (corresponding author) is supported by the 111 Project (B14019) and Shanghai Pujiang Program (14PJ1403300). This project has been carried out while S. Tindel was on sabbatical at the University of Kansas. He wishes to express his gratitude to this institution for its warm hospitality.

References

- [1] ARCONES, M.A. (1994). Limit theorems for nonlinear functionals of a stationary Gaussian sequence of vectors. *Ann. Probab.* **22** 2242–2274. [MR1331224](#)
- [2] BERAN, J. (1994). *Statistics for Long-Memory Processes. Monographs on Statistics and Applied Probability* **61**. New York: Chapman and Hall. [MR1304490](#)
- [3] BREUER, P. and MAJOR, P. (1983). Central limit theorems for nonlinear functionals of Gaussian fields. *J. Multivariate Anal.* **13** 425–441. [MR0716933](#)
- [4] BROCKWELL, P.J. and DAVIS, R.A. (1991). *Time Series: Theory and Methods*, 2nd ed. *Springer Series in Statistics*. New York: Springer. [MR1093459](#)
- [5] CARBERY, A. and WRIGHT, J. (2001). Distributional and L^q norm inequalities for polynomials over convex bodies in \mathbb{R}^n . *Math. Res. Lett.* **8** 233–248. [MR1839474](#)

- [6] CHAMBERS, D. and SLUD, E. (1989). Central limit theorems for nonlinear functionals of stationary Gaussian processes. *Probab. Theory Related Fields* **80** 323–346. [MR0976529](#)
- [7] GIRAITIS, L. and SURGAILIS, D. (1985). CLT and other limit theorems for functionals of Gaussian processes. *Z. Wahrsch. Verw. Gebiete* **70** 191–212. [MR0799146](#)
- [8] GRANGER, C.W.J. and JOYEUX, R. (1980). An introduction to long-memory time series models and fractional differencing. *J. Time Series Anal.* **1** 15–29. [MR0605572](#)
- [9] GUBNER, J.A. (2005). Theorems and fallacies in the theory of long-range-dependent processes. *IEEE Trans. Inform. Theory* **51** 1234–1239. [MR2237996](#)
- [10] HOSKING, J.R.M. (1981). Fractional differencing. *Biometrika* **68** 165–176. [MR0614953](#)
- [11] HU, Y., LU, F. and NUALART, D. (2014). Convergence of densities of some functionals of Gaussian processes. *J. Funct. Anal.* **266** 814–875. [MR3132731](#)
- [12] NOURDIN, I. and NUALART, D. (2013). Fisher information and the fourth moment theorem. Preprint. Available at [arXiv:1312.5841](#).
- [13] NOURDIN, I. and PECCATI, G. (2009). Stein’s method on Wiener chaos. *Probab. Theory Related Fields* **145** 75–118. [MR2520122](#)
- [14] NOURDIN, I. and PECCATI, G. (2012). *Normal Approximations with Malliavin Calculus: From Stein’s Method to Universality*. *Cambridge Tracts in Mathematics* **192**. Cambridge: Cambridge Univ. Press. [MR2962301](#)
- [15] NOURDIN, I., PECCATI, G. and PODOLSKIJ, M. (2011). Quantitative Breuer–Major theorems. *Stochastic Process. Appl.* **121** 793–812. [MR2770907](#)
- [16] NOURDIN, I. and POLY, G. (2013). Convergence in total variation on Wiener chaos. *Stochastic Process. Appl.* **123** 651–674. [MR3003367](#)
- [17] NUALART, D. (2006). *The Malliavin Calculus and Related Topics*, 2nd ed. *Probability and Its Applications (New York)*. Berlin: Springer. [MR2200233](#)
- [18] NUALART, D. and PECCATI, G. (2005). Central limit theorems for sequences of multiple stochastic integrals. *Ann. Probab.* **33** 177–193. [MR2118863](#)
- [19] PALMA, W. (2007). *Long-Memory Time Series: Theory and Methods*. *Wiley Series in Probability and Statistics*. Hoboken, NJ: Wiley. [MR2297359](#)

Received March 2014 and revised May 2014